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Asymptotic Analysis of Sample Average Approximation for Stochastic Optimization Problems with Joint Chance Constraints via Conditional Value at Risk and Difference of Convex Functions

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Abstract Conditional Value at Risk (CVaR) has been recently used to approximate a chance constraint. In this paper, we study the convergence of stationary points, when sample average approximation (SAA) method is applied to a CVaR approximated joint chance constrained stochastic minimization problem. Specifically, we prove under some moderate conditions that optimal solutions and stationary points, obtained from solving sample average approximated problems, converge with probability one to their true counterparts. Moreover, by exploiting the recent results on large deviation of random functions and sensitivity results for generalized equations, we derive exponential rate of convergence of stationary points. The discussion is also extended to the case, when CVaR approximation is replaced by a difference of two convex functions (DC-approximation). Some preliminary numerical test results are reported.

Keywords Joint chance constraints · CVaR · DC-approximation · Almost H-calmness · Stationary point · Exponential convergence

1 Introduction

Joint chance constrained optimization models have wide applications in communication and networks, product design, system control, statistics, and finance; see [1] for details. It is well known that joint chance constrained optimization problems are

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difficult to solve as the probabilistic function in the chance constraint is generally nonconvex, and it is difficult to obtain a closed form or to be evaluated. Various approaches have been proposed in the literature to tackle the difficulties. For example, a number of convex conservative approximation schemes, such as quadratic approximation [2], CVaR approximation [3] and Bernstein approximation [4], have been proposed for the chance constraint. CVaR approximation seems to be one of the most widely used approximation schemes since it is numerically tractable and enjoys nice features such as convexity and monotonicity. The approximation scheme was first proposed by Nemirovski and Shapiro [4] for single chance constrained minimization problem and has now been widely applied in stochastic programming.

In this paper, we consider a stochastic minimization problem with a joint chance constraint, which is approximated by CVaR. We study numerical methods for the approximated problem. The main challenge here is to handle the expected value of the underlying functions. If we are able to obtain a closed form of the mathematical expectation, then the problem becomes an ordinary nonlinear programming problem (NLP). However, in practice, it is often difficult to do so, not only because there is a max operation in the integrand, but because it requires complete information of the distribution of the underlying random variables or multidimensional integration. A well-known method to tackle this problem is sample average approximation (SAA), which is also widely known under various names such as Monte Carlo method, sample path optimization (SPO) [5], and stochastic counterpart; see [6] for a comprehensive review. The basic idea of SAA is to use sample average to approximate the expected value. For a fixed sample, the sample average approximated problem is a deterministic nonlinear problem (NLP) and, therefore, any appropriate NLP code can be applied to solve the problem.

Our focus in this paper is to analyze whether an optimal solution or a stationary point obtained from solving the SAA problem converges to their true counterpart as sample size increases, and if so at what rate. The latter is practically interesting as one would like to know how large the sample size should be in order to obtain an approximate solution within a specified precision. This kind of analysis is known as asymptotic convergence analysis in stochastic programming, and it is technically challenging in that CVaR is in the constraint rather than at the objective. The latter has been well studied; see, for instance, [7, 8].

In order to carry out the convergence analysis, particularly in relation to stationary points, we need to derive first order optimality conditions. We do so under convex and nonconvex settings and resort to Hiriart-Urruty's earlier results on Karush–Kuhn–Tucker (KKT for short) conditions in nonsmooth constrained optimization [9].

In the literature of continuous stochastic optimization, asymptotic analysis of a statistical estimator of optimal value and optimal solution is based on uniform exponential convergence of a random function which is Hölder continuous; see, for instance, [10, 11]. Xu [12] extends the results to a class of so-called H-calm functions which allow some extent of discontinuity and this is used to derive exponential convergence of stationary points in nonsmooth stochastic optimization. In a more recent development, the exponential convergence results are further extended [13] to a class of almost H-calm functions. Here, we present a detailed discussion about the advantage of almost H-calmness (see Remark 4.1 and Example 4.1) and strengthen

[13, Theorem 3.1] by weakening a boundedness condition imposed on the random function (see Theorem 8.1). We then apply the strengthened uniform convergence result to establish exponential convergence of Clarke stationary points (Theorem 4.2). We also do so for the case when the chance constraint is approximated by a DC-approximation. This strengthens the existing results by Hong et al. [14].

The rest of the paper is organized as follows. Section 2 describes the setup of the original problem, its approximation schemes, and calculus of Clarke subdifferentials of the underlying functions. Section 3 discusses optimality conditions of true problem and the SAA problem for CVaR approximated problems in both convex and nonconvex cases. Section 4 is devoted to a detailed asymptotic convergence analysis of stationary points of the SAA problem as sample size increases, and Sect. 5 presents similar convergence analysis to DC-approximated problem. Section 6 reports some preliminary numerical test results, and finally Sect. 7 concludes with some remarks.

2 Preliminaries

2.1 Notation

Throughout this paper, we use the following notation. $x^T y$ denotes the scalar product of two vectors x and y , $x \perp y$ denotes their perpendicularity, and x^T denotes its transposition. $\|\cdot\|$ denotes the Euclidean norm of a vector. We write $d(x, D) := \inf_{x' \in D} \|x - x'\|$ for the distance from point x to set D . For two sets D_1 and D_2 ,

$$\mathbb{D}(D_1, D_2) := \sup_{x \in D_1} d(x, D_2)$$

denotes the deviation of set D_1 from set D_2 . For a real valued function $h(x)$, we use $\nabla h(x)$ to denote the gradient of h at x . If $h(x)$ is vector valued, then the same notation refers to the classical Jacobian of h at x . For $x \in X$, $B(x; \rho)$ denotes a closed ρ -neighborhood of x relative to X .

Let $v : \mathbb{R}^n \rightarrow \mathbb{R}^m$ be a locally Lipschitz continuous function. Recall that *Clarke generalized derivative* of v at point x in direction d is defined as

$$v^o(x, d) := \limsup_{y \rightarrow x, t \downarrow 0} \frac{v(y + td) - v(y)}{t}.$$

v is said to be *Clarke regular* at x if the usual one sided directional derivative, denoted by $v'(x, d)$, exists for every $d \in \mathbb{R}^n$ and $v^o(x, d) = v'(x, d)$. The *Clarke generalized gradient* (also known as Clarke subdifferential) is defined as

$$\partial v(x) := \{\zeta : \zeta^T d \leq v^o(x, d)\}.$$

For a vector valued function v , it is said to be strictly differentiable at x , if v admits a strict derivative at x , an element of $\mathcal{L}(\mathbb{R}^n, \mathbb{R}^m)$ denoted $Dv(x)$, proved for each d , the following holds:

$$\lim_{y \rightarrow x, t \downarrow 0} \frac{v(y + td) - v(y)}{t} = d^T Dv(y),$$

and provided the convergence is uniform for d in compact sets, where $\mathcal{L}(\mathbb{R}^n, \mathbb{R}^m)$ denote a space of continuous linear function from \mathbb{R}^n to \mathbb{R}^m . Note that, if v is continuously differentiable, it is strictly differentiable; see [15, Chap. 2]. Note that, in this paper, we need to consider Clarke subdifferential of a random function $v(x, \xi)$ with respect to x , that is, $\partial_x v(x, \xi)$. In such a case, it is meant the Clarke subdifferential of function $v(\cdot, \xi)$ at point x for fixed ξ . It is well known that when $v(x, \xi)$ is convex w.r.t. x , the Clarke subdifferential coincides with the subdifferential in the sense of convex analysis; see [15, Proposition 2.2.7]. Note that unless the function is strictly differentiable, Clarke differential does not collapse to the classic one.

2.2 Problem Statement

Consider the following joint chance constrained minimization problem:

$$\min_{x \in \mathcal{X}} f(x) \quad \text{s.t.} \quad \text{Prob}\{c(x, \xi) \leq 0\} \geq 1 - \alpha, \quad (1)$$

where \mathcal{X} is a convex compact subset of \mathbb{R}^n , ξ is a random vector in a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ with support $\mathcal{E} \subset \mathbb{R}^k$ and α is a positive number between 0 and 1; $f: \mathbb{R}^n \rightarrow \mathbb{R}$ is a continuously differentiable function, $c(x, \xi) := \max\{c_1(x, \xi), \dots, c_m(x, \xi)\}$ and $c_i: \mathbb{R}^n \times \mathbb{R}^k \rightarrow \mathbb{R}$, $i = 1, \dots, m$, is continuously differentiable w.r.t. x for every fixed ξ .

In what follows, we focus on CVaR approximation method for problem (1). Recall that Value-at-Risk (VaR) of a random function $c(x, \xi)$ is defined as

$$\text{VaR}_{1-\alpha}(c(x, \xi)) := \min_{\eta \in \mathbb{R}} \{\eta : \text{Prob}\{c(x, \xi) \leq \eta\} \geq 1 - \alpha\}$$

and CVaR is defined as the conditional expected value of $c(x, \xi)$ exceeding VaR, that is,

$$\text{CVaR}_{1-\alpha}(c(x, \xi)) := \frac{1}{\alpha} \int_{c(x, \xi) \geq \text{VaR}_{1-\alpha}(c(x, \xi))} c(x, \xi) P(d\xi).$$

The latter can be reformulated as

$$\min_{\eta \in \mathbb{R}} \left(\eta + \frac{1}{\alpha} \mathbb{E}[(c(x, \xi) - \eta)_+] \right);$$

see [3]. It is easy to verify that $\text{CVaR}_{1-\alpha}(c(x, \xi)) \rightarrow \text{VaR}_{1-\alpha}(c(x, \xi))$ as $\alpha \downarrow 0$. In the case when $c(x, \xi)$ is convex in x for almost every ξ , $\text{CVaR}_{1-\alpha}(c(x, \xi))$ is a convex function; see [3, Theorem 2]. Observe that the chance constraint in problem (1) can be written as $\text{VaR}_{1-\alpha}(c(x, \xi)) \leq 0$ while CVaR is often regarded as an approximation of VaR. Therefore, we may consider the following approximation scheme for (1) by replacing the chance constraint with a CVaR constraint:

$$\min_{x \in \mathcal{X}} f(x) \quad \text{s.t.} \quad \text{CVaR}_{1-\alpha}(c(x, \xi)) \leq 0. \quad (2)$$

Using the reformulation of CVaR, i.e., [3], we may present problem (2) as

$$\min_{(x, \eta) \in \mathcal{X} \times \mathbb{R}} f(x) \quad \text{s.t.} \quad \eta + \frac{1}{\alpha} \mathbb{E}[(c(x, \xi) - \eta)_+] \leq 0; \quad (3)$$

see for instance [16, Sect. 1.4]. Moreover, we apply the SAA method to deal with the expect value. Let ξ^1, \dots, ξ^N be an independent and identically distributed (i.i.d.) sampling of ξ . We consider the following sample average approximation problem for problem (3):

$$\min_{(x, \eta) \in \mathcal{X} \times \mathbb{R}} f(x) \quad \text{s.t.} \quad \eta + \frac{1}{\alpha N} \sum_{j=1}^N (c(x, \xi^j) - \eta)_+ \leq 0. \quad (4)$$

We refer to (3) as the *true* problem and (4) as its SAA problem.

For the simplicity of notation, let

$$\begin{aligned} p(x) &:= \max(0, x), & g(x, \eta, \xi) &:= p(c(x, \xi) - \eta), \\ G(x, \eta) &:= \mathbb{E}[g(x, \eta, \xi)], & G_N(x, \eta) &:= \frac{1}{N} \sum_{j=1}^N g(x, \eta, \xi^j), \\ H(x, \eta) &:= \eta + \frac{1}{\alpha} G(x, \eta), & \text{and} \quad H_N(x, \eta) &:= \eta + \frac{1}{\alpha} G_N(x, \eta). \end{aligned}$$

Assumption 2.1 There exists a point $x_0 \in \mathcal{X}$ such that $\mathbb{E}[(c(x_0, \xi))_+] < \infty$. Moreover,

$$\mathbb{E}[\|\nabla_x c_i(x, \xi)\|] < \infty, \quad \text{for } i = 1, \dots, m, x \in \mathcal{X}.$$

Proposition 2.1 Suppose that Assumption 2.1 holds. Then

- (i) $H(x, \eta)$ is well defined for all $x \in \mathcal{X}$ and $\eta \in \mathbb{R}$, locally Lipschitz continuous w.r.t. x , globally Lipschitz continuous w.r.t. η and

$$\partial_{(x, \eta)} H(x, \eta) \subset \begin{pmatrix} 0 \\ 1 \end{pmatrix} + \mathbb{E}[\partial_{(x, \eta)} g(x, \eta, \xi)], \quad (5)$$

where

$$\partial_{(x, \eta)} g(x, \eta, \xi) \subset \partial p(c(x, \xi) - \eta) (\partial_x c(x, \xi), -1)^T, \quad (6)$$

and $\partial_x c(x, \xi) = \text{conv}\{\nabla_x c_i(x, \xi), i \in i(x)\}$, $i(x) := \{i : c(x, \xi) = c_i(x, \xi)\}$,

$$\partial p(\eta - c(x, \xi)) = \begin{cases} 1, & \text{if } c(x, \xi) - \eta > 0, \\ [0, 1], & \text{if } c(x, \xi) - \eta = 0, \\ 0, & \text{if } c(x, \xi) - \eta < 0, \end{cases}$$

and the expected value of the Clarke subdifferential of the random function is in the sense of Aumann's integral [17], here and later on, "conv" denotes convex hull of a set;

- (ii) if $c(x, \xi)$ is convex or strictly differentiable w.r.t. x for all $\xi \in \Xi$, then $H(x, \eta)$ is Clarke regular w.r.t. x and η and the equality in (5) and (6) holds.

Proof Part (i). Verification of the well definedness and Lipschitzness is elementary given the fact that $p(c(x, \xi) - \eta)$ is a composition of the max function $p(\cdot)$ and $c(x, \xi) - \eta$. In what follows, we show inclusions (5) and (6). By the definition of Clarke subdifferential,

$$\partial_{(x, \eta)}(c(x, \xi) - \eta) = (\partial_x c(x, \xi), -1)^T.$$

Since $c_i(x, \xi)$ is continuously differentiable, by [15, Proposition 2.3.12],

$$\partial_x c(x, \xi) = \text{conv}\{\nabla_x c_i(x, \xi), i \in i(x)\},$$

where $i(x) := \{i : c(x, \xi) = c_i(x, \xi)\}$. Through the chain rule ([15, Theorem 2.3.9]), we obtain

$$\partial_{(x, \eta)} g(x, \eta, \xi) \subset \partial p(c(x, \xi) - \eta)(\partial_x c(x, \xi), -1)^T.$$

Furthermore, it is easy to verify that the term at the right-hand side of the formula above is bounded by $\max_{i=1}^m \{\|\nabla_x c_i(x, \xi)\|\} + 1$ which is integrably bounded under Assumption 2.1. Following a discussion by Artstein and Vitale [18], both $\mathbb{E}[\partial_x c(x, \xi)]$ and $\mathbb{E}[\partial p(c(x, \xi) - \eta)]$ are well defined. Further, by [12, Theorem 2.1],

$$\partial_{(x, \eta)} G(x, \eta) \subset \partial \mathbb{E}[p(c(x, \xi) - \eta)(\partial_x c(x, \xi), -1)^T]$$

which implies (5).

Part (ii). In the case when $c(x, \xi)$ is convex or strictly differentiable w.r.t. x for all $\xi \in \Xi$, $g(x, \eta, \xi)$ is Clarke regular and equality in (6) holds. Subsequently, we have

$$\partial_{(x, \eta)} G(x, \eta) = \mathbb{E}[\partial p(c(x, \xi) - \eta)(\partial_x c(x, \xi), -1)^T]$$

and hence equality in (5) holds. \square

3 Optimality Conditions

In this section, we discuss optimality conditions and first order necessary conditions of the true problems (3) and its sample average approximated counterpart. This is to pave the way for the asymptotic convergence analysis of stationary points of problem (4) as sample size increases.

3.1 Optimality Conditions of the True Problem

Let us start with true problems. A widely used condition for deriving optimality conditions of a constrained convex program is Slater's constraint qualification.

Assumption 3.1 Problem (3) satisfies the *Slater's constraint qualification*, that is, there exists a point $(x_0, \eta_0) \in \mathcal{X} \times \mathbb{R}$ such that $H(x_0, \eta_0) < 0$.

Let $\lambda \geq 0$ be a number and define the Lagrange function of problem (3):

$$\mathcal{L}(x, \eta, \lambda) := f(x) + \lambda H(x, \eta).$$

Proposition 3.1 Assume that f and $c_i, i = 1, \dots, m$, are convex w.r.t. x . Let (x^*, η^*) be an optimal solution of (3). Under Assumptions 2.1–3.1, there exists a number $\lambda^* \in \mathbb{R}_+$ such that

$$\begin{cases} (x^*, \eta^*) \in \arg \min_{(x, \eta) \in \mathcal{X} \times R} \mathcal{L}(x, \eta, \lambda^*), \\ 0 \leq -H(x^*, \eta^*) \perp \lambda^* \geq 0. \end{cases} \quad (7)$$

The set of λ^* satisfying (7) is nonempty, convex, and bounded, and is the same for any optimal solution of the problem.

Proof Since f and $c_i, i = 1, \dots, m$ are convex functions, problem (3) is convex. By Assumption 3.1 and Bonnans and Shapiro [19, Proposition 2.106], Robinson's constraint qualification holds. The conclusion then follows from [19, Theorem 3.4]. \square

It is possible to characterize the optimality conditions (7) in terms of the derivatives of the underlying functions. In what follows, we do so for general case, by invoking to Hiriart-Urruty's KKT conditions for a nonsmooth problem with equality and inequality constraints [9]. Recall that the Bouligrand tangent cone to a set $X \subset \mathbb{R}^n$ at a point $x \in X$ is defined as follows:

$$\mathcal{T}_X(x) := \{u \in \mathbb{R}^n : d(x + tu, X) = o(t), t \geq 0\}.$$

The normal cone to X at x , denoted by $\mathcal{N}_X(x)$, is the polar of the tangent cone:

$$\mathcal{N}_X(x) := \{\zeta \in \mathbb{R}^n : \zeta^T u \leq 0, \forall u \in \mathcal{T}_X(x)\}$$

and $\mathcal{N}_X(x) = \emptyset$ if $x \notin X$.

Let $\Phi(x)$ be a locally Lipschitz function defined on an open subset $\mathcal{O} \subset \mathbb{R}^n$. Let $x_0 \in \mathcal{O}$ and let Q be a subset of \mathbb{R}^n such that $x_0 \in \bar{Q}$ (closure of Q). We denote by $\mathcal{V}_Q(x_0)$ the filter of neighborhoods of x_0 for the topology induced on Q . The collection $(\partial_x \Phi(x) | x \in \mathcal{O}, \mathcal{V}_Q(x_0))$ is a filtered family [20, p. 126]. For this family, we may consider the “lim sup,” which we will denote by $\partial_x^Q \Phi(x_0)$

$$\partial_x^Q \Phi(x_0) := \bigcap_{V \in \mathcal{V}_Q(x_0)} \overline{\bigcup_{x \in V} \partial_x \Phi(x)}.$$

In other words, $\partial_x^Q \Phi(x_0)$ consists of all cluster points of sequences of matrices $M_i \in \partial_x \Phi(x_i)$ as x_i converges to x_0 in Q . Obviously, $\partial_x^Q \Phi(x_0) \subset \partial_x \Phi(x_0)$; see [9]. The following constraint qualification stems from Hiriart-Urruty [9].

Definition 3.1 Problem (3) is said to satisfy the *subdifferential constraint qualification* at a feasible point (x, η) iff, for all $\zeta \in \partial_{(x, \eta)}^{S^c} H(x, \eta)$, there exists $d \in \mathbb{R}^{n+1}$ such that $\zeta^T d < -\delta$, where S denotes the feasible set of the problem (3) and S^c is complementary set of S .

In the case when $\delta = 0$, the condition is regarded as an extended *Mangasarian–Fromovitz constraint qualification* (MFCQ), see discussions at pp. 79–80 in Hiriart-Urruty [9]. Indeed, if $H(x, \eta)$ is differentiable and $\delta = 0$, then the subdifferential constraint qualification reduces to the classical MFCQ.

Theorem 3.1 *Let $(x^*, \eta^*) \in \mathcal{X} \times \mathbb{R}$ be a local optimal solution to the true problem (3). Suppose that Assumption 2.1 holds and the subdifferential constraint qualification is satisfied at (x^*, η^*) . Then there exists $\lambda^* \in \mathbb{R}_+$ such that (x^*, η^*, λ^*) satisfies the following Karush–Kuhn–Tucker (KKT for short) conditions*

$$\begin{cases} 0 \in \begin{pmatrix} \nabla f(x) \\ 0 \end{pmatrix} + \lambda \partial_{(x, \eta)} H(x, \eta) + \mathcal{N}_{\mathcal{X} \times \mathbb{R}}(x, \eta), \\ 0 \leq -H(x, \eta) \perp \lambda \geq 0, \end{cases} \quad (8)$$

which imply

$$\begin{cases} 0 \in \begin{pmatrix} \nabla f(x) \\ 0 \end{pmatrix} + \lambda \mathbb{E} \left[\begin{pmatrix} 0 \\ 1 \end{pmatrix} + \frac{1}{\alpha} \partial_{(x, \eta)} g(x, \eta, \xi) \right] + \mathcal{N}_{\mathcal{X} \times \mathbb{R}}(x, \eta), \\ 0 \leq -H(x, \eta) \perp \lambda \geq 0. \end{cases} \quad (9)$$

If $c(x, \xi)$ is strictly differentiable w.r.t. x , then (8) and (9) are equivalent.

Proof The KKT conditions (8) follow from Proposition 2.1 and [9, Theorem 4.2]. The conditions in (9) follow from (5) and (8).

In the case when $c(x, \xi)$ is strictly differentiable, one can easily use Proposition 2.1(ii) to show that KKT conditions (9) are equivalent to (8). The proof is complete. \square

Remark 3.1 In the case when $f(x)$ and $c(x, \xi)$ are convex w.r.t. x , the subdifferentiable constraint qualification in Theorem 3.1 can be weakened to Slater's constraint qualification, in which case, (8) and (9) are equivalent without strict differentiability. The claim follows from Proposition 3.1.

The KKT conditions (9) are weaker than those of (8) due to (5) in general; see Xu [12] for a detailed discussion on this.

A tuple $(x, \eta, \lambda) \in \mathcal{X} \times \mathbb{R} \times \mathbb{R}_+$ satisfying (8) is called a KKT pair of problem (3), (x, η) a *Clarke stationary point* and λ the corresponding Lagrange multiplier. Similarly, a tuple $(x, \eta, \lambda) \in \mathcal{X} \times \mathbb{R} \times \mathbb{R}_+$ satisfying (9) is called a weak KKT pair of problem (3), (x, η) a *weak Clarke stationary point* and λ the corresponding Lagrange multiplier.

3.2 Optimality Conditions of SAA Problem

We now move on to discuss the optimality conditions of SAA problem (4). We need the following technical results.

Proposition 3.2 *Let Assumption 2.1 holds. Let $Z =: \mathcal{X} \times \mathbb{R}$ and \mathcal{Z} be a compact subset of Z . Then*

- (i) *with probability one (w.p.1) $\frac{1}{N} \sum_{j=1}^N c(x, \xi^j)$ converge respectively to $\mathbb{E}[c(x, \xi)]$ uniformly over any compact subset of \mathcal{X} as $N \rightarrow \infty$;*

- (ii) w.p.1 $\frac{1}{N} \sum_{i=1}^N g(x, \eta, \xi)$ converges respectively to $G(x, \eta)$ uniformly on \mathcal{Z} as $N \rightarrow \infty$;
- (iii) if, in addition, Assumption 3.1 holds, then the SAA problem (4) satisfies the Slater's constraint qualification w.p.1 for N sufficiently large, that is, there exists a point $(x, \eta) \in \mathcal{X} \times \mathbb{R}$ such that $H_N(x, \eta) < 0$ w.p.1 for N sufficiently large.

Proof Part (i) and Part (ii) follow from [6, Proposition 7]. Part (iii) is straightforward from parts (i)–(ii) and the definition of Slater's constraint qualification. \square

From Proposition 3.2, we know that for N sufficiently large, the SAA problem (4) satisfies the Slater's constraint qualification w.p.1. Consequently, we have the following optimality conditions for problem (4).

Theorem 3.2 Assume that f and $c_i, i = 1, \dots, m$, are convex w.r.t. x and Assumptions 2.1–3.1 hold. If (x_N, η_N) is an optimal solution of the problem (4), then there exists a number $\lambda_N \in \mathbb{R}_+$ such that w.p.1

$$\begin{cases} (x_N, \eta_N) \in \arg \min_{(x, \eta) \in \mathcal{X} \times \mathbb{R}} \mathcal{L}_N(x, \eta, \lambda_N), \\ 0 \leq -H_N(x_N, \eta_N) \perp \lambda_N \geq 0. \end{cases} \quad (10)$$

Moreover, w.p.1 (x_N, η_N, λ_N) satisfies the KKT conditions

$$\begin{cases} 0 \in \left(\nabla f_0(x) \right) + \lambda \partial_{(x, \eta)} H_N(x, \eta) + \mathcal{N}_{\mathcal{X} \times \mathbb{R}}(x, \eta), \\ 0 \leq -H_N(x, \eta) \perp \lambda \geq 0, \end{cases} \quad (11)$$

which imply

$$\begin{cases} 0 \in \left(\nabla f_0(x) \right) + \lambda \left(\begin{pmatrix} 0 \\ 1 \end{pmatrix} + \frac{1}{\alpha N} \sum_{i=1}^N \partial_{(x, \eta)} g(x, \eta, \xi) \right) + \mathcal{N}_{\mathcal{X} \times \mathbb{R}}(x, \eta), \\ 0 \leq -H_N(x, \eta) \perp \lambda \geq 0. \end{cases} \quad (12)$$

and the set of λ_N satisfying (10) is nonempty, convex, and bounded, and is the same for any optimal solution of the problem.

Proof By Proposition 3.2, the Slater's condition holds for problem (4). Similar to Proposition 3.1, (10) holds, which implies (11). Since $c(x, \xi)$ is convex, by Proposition 2.1(ii), (12) and (11) are equivalent. The proof is complete. \square

In what follows, we derive the KKT conditions for nonconvex case. We need some additional conditions.

Definition 3.2 Problem (3) is said to satisfy the *strong subdifferential constraint qualification* at a feasible point (x, η) iff there exist a positive number δ and $d \in \mathbb{R}^{n+1}$ such that $\zeta^T d < -\delta$ for all $\zeta \in \partial_{(x, \eta)} H(x, \eta)$.

Note that, since $\partial_{(x, \eta)}^{S^c} H(x, \eta) \subset \partial_{(x, \eta)} H(x, \eta)$, the strong subdifferential constraint qualification implies the subdifferential constraint qualification. The terminology was introduced by Dentcheva and Ruszczyński in [21] for a class of semiinfinite

optimization problems. Through the separation theorem in convex analysis, it is easy to observe that the constraint qualification is equivalent to $0 \notin \partial H(x, \eta)$.

Theorem 3.3 *Let $(x_N, \eta_N) \in \mathcal{X} \times \mathbb{R}$ be a local optimal solution to the sample average approximation problem (4). Let \hat{Z} denote the subset of points in $\mathcal{X} \times \mathbb{R}$ such that*

$$\lim_{N \rightarrow \infty} d((x_N, \eta_N), \hat{Z}) \rightarrow 0$$

w.p.1. Assume that \hat{Z} is bounded, Assumption 2.1 holds and problem (3) satisfies the strong subdifferential constraint qualification at every point $(x, \eta) \in \hat{Z}$. Then w.p.1 there exists $\lambda_N \in R_+$ such that (x_N, η_N, λ_N) satisfies the KKT conditions (11) which imply (12). Moreover, if $c(x, \xi)$ is strictly differentiable w.r.t. x for all $\xi \in \Xi$, then KKT conditions (11) are equivalent to that of (12).

Proof Under the strong subdifferential constraint qualification at (x, η) , there exist a constant $\delta > 0$ and a vector $u \neq 0$ (which depends on (x, η)) such that

$$\sup_{\zeta \in \partial_{(x, \eta)} H(x, \eta)} \zeta^T u = H^o(x, \eta; u) \leq -\delta.$$

In what follows, we show that

$$\sup_{\zeta \in \partial_{(x, \eta)} H_N(x_N, \eta_N)} \zeta^T u \leq -\delta/2 \quad (13)$$

w.p.1 for N sufficiently large, i.e., problem (4) satisfies the subdifferential constraint qualification. Let $(x, \eta) \in \hat{Z}$,

$$\partial_{(x, \eta)}^\epsilon H(x, \eta) := \bigcup_{(x', \eta') \in (x, \eta) + \epsilon \mathcal{B}} \partial_{(x', \eta')} H(x', \eta')$$

and

$$H_\epsilon^o(x, \eta; u) := \sup_{\zeta \in \partial_{(x, \eta)}^\epsilon H(x, \eta)} \zeta^T u.$$

The outer semicontinuity of $\partial_{(x, \eta)} H(\cdot, \cdot)$ allows us to find a sufficiently small ϵ (depending on (x, η)) such that

$$H_\epsilon^o(x', \eta'; u) \leq -\frac{3}{4}\delta$$

for all $(x', \eta') \in B(x, \eta; \rho)$. Let \tilde{Z} denote the closed ρ neighborhood of \hat{Z} relative to $\mathcal{X} \times \mathbb{R}$. Applying [22, Lemma 2.1] to \tilde{Z} , we can find a positive number \hat{N} such that

$$H_N^o(x', \eta'; u) - H_\epsilon^o(x', \eta'; u) \leq \frac{\delta}{2}, \quad \forall (x', \eta') \in \tilde{Z}$$

w.p.1 for $N \geq \hat{N}$. Using this inequality, we have

$$\sup_{\zeta \in \partial_{(x, \eta)} H_N(x_N, \eta_N)} \zeta^T u = H_N^o(x_N, \eta_N; u) \leq H_\epsilon^o(x_N, \eta_N; u) + \frac{\delta}{2} \leq -\frac{\delta}{4}$$

w.p.1 as long as $(x_N, \eta_N) \in B(x, \eta; \rho)$. That means for all $\zeta \in \partial_{(x, \eta)} H_N(x_N, \eta_N)$, w.p.1 $\zeta^T u \leq -\frac{\delta}{4}$. It implies that w.p.1 problem (4) satisfies the strong subdifferential constraint qualification at z_N for N sufficiently large. The rest of the proof is similar to Theorem 3.1. \square

A triad $(x, \eta, \lambda) \in \mathcal{X} \times \mathbb{R} \times \mathbb{R}_+$ satisfying (12) is called a weak KKT pair of problem (4) and a tuple $(x, \eta, \lambda) \in \mathcal{X} \times \mathbb{R} \times \mathbb{R}_+$ satisfying (11) is called a KKT pair of problem (4).

We make a blanket assumption that throughout the rest of the paper the conditions of Theorems 3.2 or 3.3 hold.

4 Convergence Analysis

In this section, we discuss convergence of SAA problem (4) as sample size N increases. Differing from many asymptotic analysis in the literature, our focus here is on the convergence of stationary points/KKT pair of the SAA problem in that a local or global optimal solution is also a stationary point. The analysis is practically useful in that: (a) when the problem is nonconvex, it is often difficult to obtain a global or even a local optimal solution, a stationary point might provide some information on local optimality; (b) CVaR approximation problem (2) has potential applications in finance and engineering. Our analysis is divided into two parts: almost sure convergence and exponential convergence. The former is to examine the asymptotic consistency of stationary points obtained from solving the SAA problem and the latter is to investigate the rate of convergence through large deviation theorem. Throughout this section, we assume that the probability space (Ω, \mathcal{F}, P) is nonatomic.

4.1 Almost Sure Convergence

Consider the SAA problem (4). Assume that, for each given sampling, we solve the problem and obtain a stationary point (x_N, η_N) which satisfies (12). We investigate the convergence of (x_N, η_N) as N increases.

Assumption 4.1 Let $Z := \mathcal{X} \times \mathbb{R}$. There exists a compact subset $\mathcal{Z} \times \Lambda \subset Z \times \mathbb{R}_+$ and a positive number N_0 such that w.p.1 problem (4) has a KKT pair $(x_N, \eta_N, \lambda_N) \in \mathcal{Z} \times \Lambda$ for $N \geq N_0$.

This assumption is standard; see, for instance, Ralph and Xu [23].

Recall that for a set D , the support function of D is defined as

$$\sigma(D, u) := \sup_{d \in D} d^T u.$$

Let D_1, D_2 be two convex and compact subsets of \mathbb{R}^m . Let $\sigma(D_1, u)$ and $\sigma(D_2, u)$ denote the support functions of D_1 and D_2 , respectively. Then

$$\mathbb{D}(D_1, D_2) = \max_{\|u\| \leq 1} (\sigma(D_1, u) - \sigma(D_2, u)). \quad (14)$$

The above relationship is known as Hörmander's formula; see [24, Theorem II-18].

Theorem 4.1 *Let $\{(x_N, \eta_N, \lambda_N)\}$ be a sequence of KKT pairs satisfying (12) and (x^*, η^*, λ^*) be a cluster point. Suppose Assumption 4.1 holds. Then w.p.1 (x^*, η^*, λ^*) is a weak KKT pair of the true problem (3). Moreover, if $c(x, \xi)$ is convex or strictly differentiable w.r.t. x , then w.p.1 (x^*, η^*, λ^*) is a KKT pair.*

Proof Assume without any loss of generality that (x_N, η_N, λ_N) converges to (x^*, η^*, λ^*) w.p.1 as $N \rightarrow \infty$. Since $0 \leq -H(x, \eta) \perp \lambda \geq 0$ is equivalent to $\max(\lambda, -H(x, \eta)) = 0$, in view of weak KKT conditions (9) and (12), it suffices to show that w.p.1

$$\begin{aligned} & \lim_{N \rightarrow \infty} \left\{ \begin{pmatrix} \nabla f(x_N) \\ 0 \end{pmatrix} + \left[\begin{pmatrix} 0 \\ 1 \end{pmatrix} + \frac{1}{\alpha N} \sum_{i=1}^N \partial_{(x, \eta)} g(x_N, \eta_N, \xi^i) \right] \lambda_N + \mathcal{N}_Z(x_N, \eta_N) \right\} \\ & \subset \begin{pmatrix} \nabla f(x^*) \\ 0 \end{pmatrix} + \mathbb{E} \left[\begin{pmatrix} 0 \\ 1 \end{pmatrix} + \frac{1}{\alpha} \partial_{(x, \eta)} g(x^*, \eta^*, \xi) \right] \lambda^* + \mathcal{N}_Z(x^*, \eta^*) \end{aligned} \quad (15)$$

and

$$\lim_{N \rightarrow \infty} \max(-H_N(x_N, \eta_N), \lambda_N) = \max(-H(x^*, \eta^*), \lambda^*). \quad (16)$$

Since $\partial_{(x, \eta)} g(x, \eta, \xi)$ is outer semicontinuous w.r.t. (x, η) for almost every ξ and integrably bounded, by [25, Theorem 4],

$$\sup_{(x, \eta) \in \mathcal{Z}} \mathbb{D} \left(\frac{1}{N} \sum_{i=1}^N \partial_{(x, \eta)} g(x, \eta, \xi^i), \mathbb{E}[\partial_{(x, \eta)} g(x, \eta, \xi)] \right) \rightarrow 0$$

w.p.1 as $N \rightarrow \infty$. The uniform convergence and the continuous differentiability of $f(x)$ imply (15).

On the other hand, under Assumption 2.1, it follows from Proposition 3.2 that $H_N(x, \eta)$ converges uniformly to $H(x, \eta)$ over \mathcal{Z} , which implies (16).

When $c(x, \xi)$ is convex or strictly differentiable, H is Clarke regular, hence the KKT systems (8) and (9) are equivalent. \square

4.2 Exponential Rate of Convergence

We now move on to investigate the rate of convergence to strengthen the results established in Theorem 4.1 through large deviation theorems. Results presented as such are also known as exponential convergence. While the a.s. convergence is derived fairly easily through uniform law of large numbers for random set-valued mappings, exponential convergence is far more challenging. The main technical difficulty is to

estimate the rate of uniform upper semiconvergence of sample average random set-valued mappings and we propose to deal with it by exploiting a new large deviation result in [13].

Definition 4.1 (Almost H-calmness) Let $\phi : \mathbb{R}^n \times \mathcal{E} \rightarrow \mathbb{R}$ be a real valued function and $\xi : \Omega \rightarrow \mathcal{E} \subset \mathbb{R}^k$ be a random vector defined on probability space (Ω, \mathcal{F}, P) . Let $\mathcal{X} \subset \mathbb{R}^n$ be a closed subset of \mathbb{R}^n and $x \in \mathcal{X}$ be fixed. ϕ is said to be

- (a) *almost H-calm at x from above* with modulus $\kappa_x(\xi)$ and order γ_x iff, for any $\epsilon > 0$, there exist an integrable function $\kappa_x : \mathcal{E} \rightarrow \mathbb{R}_+$, positive numbers $\gamma_x, \delta_x(\epsilon), K$ and an open set $\mathcal{E}_x(\epsilon) \subset \mathcal{E}$ such that

$$\text{Prob}(\xi \in \mathcal{E}_x(\epsilon)) \leq K\epsilon \quad (17)$$

and

$$\phi(x', \xi) - \phi(x, \xi) \leq \kappa_x(\xi) \|x' - x\|^{\gamma_x} \quad (18)$$

for all $\xi \in \mathcal{E} \setminus \mathcal{E}_x(\epsilon)$ and all $x' \in B(x, \delta_x) \cap \mathcal{X}$;

- (b) *almost H-calm at x from below* with modulus $\kappa_x(\xi)$ and order γ_x iff, for any $\epsilon > 0$, there exist an integrable function $\kappa_x : \mathcal{E} \rightarrow \mathbb{R}_+$, positive numbers $\gamma_x, \delta_x(\epsilon), K$ and an open set $\mathcal{E}_x(\epsilon) \subset \mathcal{E}$ such that

$$\text{Prob}(\xi \in \mathcal{E}_x(\epsilon)) \leq K\epsilon \quad (19)$$

and

$$\phi(x', \xi) - \phi(x, \xi) \geq -\kappa_x(\xi) \|x' - x\|^{\gamma_x} \quad (20)$$

for all $\xi \in \mathcal{E} \setminus \mathcal{E}_x(\epsilon)$ and all $x' \in B(x, \delta_x) \cap \mathcal{X}$;

- (c) *almost H-calm at x* with modulus $\kappa_x(\xi)$ and order γ_x iff, for any $\epsilon > 0$, there exist an integrable function $\kappa_x : \mathcal{E} \rightarrow \mathbb{R}_+$, positive numbers $\gamma_x, \delta_x(\epsilon), K$ and an open set $\mathcal{E}_x(\epsilon) \subset \mathcal{E}$ such that

$$\text{Prob}(\xi \in \mathcal{E}_x(\epsilon)) \leq K\epsilon \quad (21)$$

and

$$|\phi(x', \xi) - \phi(x, \xi)| \leq \kappa_x(\xi) \|x' - x\|^{\gamma_x} \quad (22)$$

for all $\xi \in \mathcal{E} \setminus \mathcal{E}_x(\epsilon)$ and all $x' \in B(x, \delta_x) \cap \mathcal{X}$.

Remark 4.1 The concept of almost H-calmness is recently proposed by Sun and Xu [13] to derive a uniform large deviation theorem for a class of discontinuous random functions ([13, Theorem 3.1]) where the underlying random variable satisfies a continuous distribution. It is closely related to the following calmness condition suggested by a referee:

- There exist an integrable function $\kappa_x : \mathcal{E} \rightarrow \mathbb{R}_+$, positive numbers γ_x, δ_x and a measurable subset $\mathcal{E}_x \subset \mathcal{E}$ such that

$$\text{Prob}(\xi \in \mathcal{E}_x) = 0 \quad (23)$$

and

$$|\phi(x', \xi) - \phi(x, \xi)| \leq \kappa_x(\xi) \|x' - x\|^{\gamma_x} \quad (24)$$

for all $\xi \in \mathcal{E} \setminus \mathcal{E}_x$ and all $x' \in B(x, \delta_x) \cap \mathcal{X}$.

Conditions (23)–(24) require H-calmness (24) to hold for *almost every* $\xi \in \mathcal{E}$ and the two conditions may be regarded as a limiting case of almost H-calmness ($\epsilon \downarrow 0_+$). Let us call the resulting calmness as *limiting almost H-calmness*. In the case when $\mathcal{E}_x = \emptyset$, it reduces to the H-calmness of [12, Definition 2.3], which requires H-calmness condition (24) to hold for *every* $\xi \in \mathcal{E}$.

Let μ denote the Lebesgue measure relative to \mathcal{E} . The limiting almost calmness conditions imply $\mu(\mathcal{E}_x) = 0$. Therefore, for any $\epsilon > 0$, there exists any open subset $\mathcal{E}_x^\epsilon \subset \mathcal{E}$ such that $\mu(\mathcal{E}_x^\epsilon) < \epsilon$. This means limiting almost H-calmness implies almost H-calmness. The reverse assertion may not be true.

The example below shows that an almost H-calm random function does not necessarily satisfy the limiting almost H-calmness condition and the necessity of almost H-calmness.

Example 4.1 Consider random function

$$\varphi(x, \xi) := \begin{cases} \frac{1}{\sqrt{|x-\xi|}}, & \text{for } x \neq \xi, \\ \infty, & \text{for } x = \xi, \end{cases}$$

where ξ is a random variable satisfying uniform distribution over $[0, 1]$. For every fixed ξ , function $\varphi(\cdot, \xi)$ is calm at any point $x \in [0, 1]$ except at point $x = \xi$ because it is locally continuously differentiable. However, this function does not satisfy limiting almost H-calmness in the sense that there does not exist positive numbers δ , γ and positive measurable function $\kappa(\xi)$ such that

$$|\varphi(x', \xi') - \varphi(x, \xi')| \leq \kappa(\xi') |x' - x|^\gamma \quad (25)$$

for all $\xi' \in \mathcal{E} \setminus \{x\}$ and $x' \in (x - \delta, x + \delta) \cap [0, 1]$. On the other hand, it is easy to verify that $\varphi(x, \xi)$ is almost H-calm at any point in $[0, 1]$. Indeed, for every $\epsilon > 0$, Let $\delta := \epsilon/2$, $\mathcal{E}_x(\epsilon) := (x - \frac{\epsilon}{2}, x + \frac{\epsilon}{2}) \cap [0, 1]$ and

$$\kappa_x(\xi) := \begin{cases} M, & \text{for } \xi \in \mathcal{E}_x \cap [0, 1], \\ \frac{1}{2\sqrt{|x-\xi|^3}}, & \text{otherwise,} \end{cases}$$

where M is any positive constant. Then it is easy to verify that

$$|\varphi(x', \xi') - \varphi(x, \xi)| \leq \kappa_x(\xi') |x' - x| \quad (26)$$

for all $\xi' \in [0, 1] \setminus \mathcal{E}_x(\epsilon)$ and $x' \in (x - \frac{1}{2}\epsilon, x + \frac{1}{2}\epsilon)$.

Note that the exponential convergence result based on almost H-calmness, [13, Theorem 3.1], requires $\varphi(x, \xi)$ to be bounded. However, it is easy to observe that

[13, Theorem 3.1] holds when $\varphi(x, \xi)$ is integrably bounded. See Theorem 8.1 in the Appendix.

In some cases, even the random function satisfies the limiting almost H-calmness condition, it is still unclear whether we are able to derive the exponential convergence through a generalization of Theorem 3.1 in [12]. To see this, consider function $\phi(x, \xi) = \max(x, \xi)$ defined on $[0, 1] \times [0, 1]$, where ξ is a random variable satisfying uniform distribution over $[0, 1]$. This function is Lipschitz continuous w.r.t. x , therefore, the Clarke generalized derivative exists. For the fixed direction $d = 1$, the derivative function can be written as

$$\phi_x^o(x, \xi; d) = \begin{cases} 1, & \text{for } x \geq \xi, \\ 0, & \text{for } x < \xi. \end{cases}$$

It is continuous w.r.t. (x, ξ) except at line $x = \xi$. In what follows, we show that it is limiting H-calm at every point in the interval $[0, 1]$. Let $x \in [0, 1]$. There exist positive constants $\gamma_x < 1$ and integrable function $\kappa_x(\cdot) := 1/\sqrt{1-x}^{\gamma_x}$ such that

$$|\phi_x^o(x', \xi'; 1) - \phi_x^o(x, \xi'; 1)| \leq \kappa_x(\xi')|x' - x|$$

for every $\xi' \in [0, 1] \setminus \{x\}$ and $x' \in [0, 1] \setminus \{x\}$. However, the moment generating function of $\kappa_x(\xi)$ defined as such does not seem to exist, which means we cannot apply [12, Theorem 3.1] to derive uniform exponential convergence for sample average approximation of the derivative function $\phi_x^o(x, \xi; 1)$ over any compact subset of $[0, 1]$. On the other hand, similar to the analysis for the random function in the preceding discussions of this example, we can show that this derivative function $\phi_x^o(x, \xi; 1)$ is indeed almost H-calm, and hence we can derive the uniform exponential convergence for the sample average of the function through Theorem 8.1. We omit the details. The discussions above effectively demonstrate the advantage of almost H-calmness over H-calmness or limiting H-calmness.

Since the subdifferential of max function take on a key role in our exponential convergence analysis and both H-calmness and limiting almost H-calmness seem too strong, we will use Theorem 8.1 which bases on almost H-calmness to derive exponential convergence of stationary points obtained from the SAA problem (4). To this end, we need a couple of technical results.

Lemma 4.1 *Let $X \in \mathbb{R}^n$ be a compact set and $x \in X$. Let $K_1(x, \xi)$ and $K_2(x, \xi)$ be continuously differentiable w.r.t. (x, ξ) and $K(x, \xi) := \max\{K_i(x, \xi), i = 1, 2\}$. Define*

$$\Xi_K(x) := \{\xi : K(x, \xi) = K_1(x, \xi) = K_2(x, \xi), \xi \in \Xi\}.$$

Assume that $\Xi_K(x)$ is compact and

$$\nabla_{\xi}(K_1(x, \xi) - K_2(x, \xi)) \neq 0, \quad \forall \xi \in \Xi_K(x).$$

Then $\mu(\Xi_K(x)) = 0$, where μ denotes the Lebesgue measure relative to Ξ . Moreover, for any $\epsilon > 0$ and any fixed $x \in X$, there exists an open set $\Xi_K^{\epsilon}(x)$ (depending on x and ϵ) such that $\Xi_K(x) \subset \Xi_K^{\epsilon}(x)$ and $\mu(\Xi_K^{\epsilon}(x) \cap \Xi) \leq \epsilon$.

Proof The conclusion follows in a straightforward way from [13, Lemma 4.1]. \square

Let $g_i(x, \eta, \xi) := c_i(x, \xi) - \eta$, for $i = 1, \dots, m$, and $g_{m+1}(x, \eta, \xi) := 0$. Then $g(x, \eta, \xi) = \max_{i=1}^{m+1} \{g_i(x, \eta, \xi)\}$. For any $i, j \in \{1, \dots, m+1\}$, $i \neq j$, define

$$\mathcal{E}_{i,j}(x, \eta) := \{\xi \in \mathcal{E} : g(x, \eta, \xi) = g_i(x, \eta, \xi) = g_j(x, \eta, \xi)\}$$

and

$$\mathcal{E}(x, \eta) := \bigcup_{i,j \in \{1, \dots, m+1\}} \mathcal{E}_{i,j}(x, \eta).$$

Obviously, $\mathcal{E}(x, \eta)$ consists of the set of $\xi \in \mathcal{E}$ such that $g(\cdot, \cdot, \xi)$ is not differentiable at (x, η) .

Proposition 4.1 *Let $\mathcal{Z} \subset \mathcal{X} \times \mathbb{R}$ be a compact set and $(x, \eta) \in \mathcal{Z}$ and ξ be a continuous random variable. Assume: (a) $c_i(x, \xi)$ is continuously differentiable w.r.t. (x, ξ) and twice continuously differentiable w.r.t. x for almost every $\xi \in \mathcal{E}$; (b) there exists an integrable function $\kappa : \mathcal{E} \rightarrow \mathbb{R}$ such that $\nabla_x c_i(\cdot, \xi)$ is Lipschitz continuous with modulus $\kappa(\xi)$ for every $\xi \in \mathcal{E}$ and $\mathbb{E}[\kappa(\xi)] < \infty$; (c) $\mathcal{E}(x, \eta)$ and $\mathcal{E}_{i,j}(x, \eta)$ are compact and*

$$\nabla_{\xi}(g_i(x, \eta, \xi) - g_j(x, \eta, \xi)) \neq 0, \quad \forall \xi \in \mathcal{E}_{i,j}(x, \eta),$$

holds for all $i, j \in \{1, \dots, m+1\}$, $i \neq j$. Then

- (i) $\mathbb{E}[g_{(x,\eta)}^o(x, \eta, \xi; u)]$ is a continuous function w.r.t. (x, η, u) ;
- (ii) if, in addition, \mathcal{E} is compact, then $g_{(x,\eta)}^o(x, \eta, \xi; u)$ is almost H -calm w.r.t. (x, η, u) with modulus $\kappa(\xi)$ and order 1 on \mathcal{Z} .

Proof Part (i). This is a well-known result. Note that condition (a) implies the twice continuous differentiability of $g_i(x, \eta, \xi)$ w.r.t. x and η and condition (b) implies locally Lipschitz continuity of $\nabla_{(x,\eta)} g_i(\cdot, \cdot, \xi)$ with modulus $\kappa(\xi)$ for every $\xi \in \mathcal{E}$. Indeed, in this case $\mathbb{E}[g(x, \eta, \xi)]$ is continuously differentiable; see, for instance, [26, Theorem 1].

Part (ii). Let $\epsilon > 0$ and $(\bar{x}, \bar{\eta}) \in \mathcal{Z}$ be fixed. Under condition (c), it follows by Lemma 4.1 that there exists an open subset $\mathcal{E}_{i,j}^{\epsilon}(\bar{x}, \bar{\eta})$ such that $\mathcal{E} \cap \mathcal{E}_{i,j}(\bar{x}, \bar{\eta}) \subset \mathcal{E}_{i,j}^{\epsilon}(\bar{x}, \bar{\eta})$ and $\mu(\mathcal{E}_{i,j}^{\epsilon}(\bar{x}, \bar{\eta})) \leq \epsilon$ for all $i, j \in \{1, \dots, m+1\}$, $i \neq j$. Let $\mathcal{E}^{\epsilon}(\bar{x}, \bar{\eta}) := \bigcup_{i,j} \mathcal{E}_{i,j}^{\epsilon}(\bar{x}, \bar{\eta})$. Since $\mathcal{E}(\bar{x}, \bar{\eta}) = \bigcup_{i,j} \mathcal{E}_{i,j}(\bar{x}, \bar{\eta})$, we have

$$\mathcal{E}(\bar{x}, \bar{\eta}) \subset \mathcal{E}^{\epsilon}(\bar{x}, \bar{\eta}) \quad \text{and} \quad \mu(\mathcal{E}^{\epsilon}(\bar{x}, \bar{\eta})) \leq \binom{m}{2} \epsilon.$$

Let $\bar{\xi} \notin \mathcal{E}^{\epsilon}(\bar{x}, \bar{\eta})$. Then there exists only a single $i \in \{1, \dots, m+1\}$ such that $g(\bar{x}, \bar{\eta}, \bar{\xi}) = g_i(\bar{x}, \bar{\eta}, \bar{\xi})$. We can find a δ -neighborhood of $(\bar{x}, \bar{\eta}, \bar{\xi})$ (depending on $\bar{x}, \bar{\eta}$ and $\bar{\xi}$), denoted by $B((\bar{x}, \bar{\eta}, \bar{\xi}), \delta_{\bar{x}, \bar{\eta}, \bar{\xi}})$ relative to $\mathcal{Z} \times \mathcal{E}$, such that for all $(x, \eta, \xi) \in B((\bar{x}, \bar{\eta}, \bar{\xi}), \delta_{\bar{x}, \bar{\eta}, \bar{\xi}})$, $g(x, \eta, \xi) = g_i(x, \eta, \xi)$. Let \bar{u} be any direction such

that $\|\bar{u}\| \leq 1$, since $g_i(x, \eta, \xi)$ is continuously w.r.t. (x, η) differentiable in a neighborhood of $(\bar{x}, \bar{\eta}, \bar{\xi})$,

$$g_{(x,\eta)}^o(x, \eta, \xi; u) = \nabla_{(x,\eta)} g_i(x, \eta, \xi)^T u.$$

Under condition (b) and the compactness of \mathcal{E} , $\nabla_{(x,\eta)} g_i(\bar{x}, \bar{\eta}, \xi)$ is bounded and there exists integrable function $\bar{\kappa}$ such that

$$\begin{aligned} & |g_{(x,\eta)}^o(x, \eta, \xi; u) - g_{(x,\eta)}^o(\bar{x}, \bar{\eta}, \xi; \bar{u})| \\ &= |\nabla_{(x,\eta)} g_i(x, \eta, \xi)^T u - \nabla_{(x,\eta)} g_i(\bar{x}, \bar{\eta}, \xi)^T \bar{u}| \\ &\leq |\nabla_{(x,\eta)} g_i(x, \eta, \xi)^T u - \nabla_{(x,\eta)} g_i(x, \eta, \xi)^T \bar{u}| \\ &\quad + |\nabla_{(x,\eta)} g_i(\bar{x}, \bar{\eta}, \xi)^T \bar{u} - \nabla_{(x,\eta)} g_i(\bar{x}, \bar{\eta}, \xi)^T \bar{u}| \\ &\leq \max_{(x,\eta)} \|\nabla_{(x,\eta)} g_i(\bar{x}, \bar{\eta}, \xi)\| \|u - \bar{u}\| + \kappa(\xi) \|\bar{u}\| \|(x, \eta) - (\bar{x}, \bar{\eta})\| \\ &\leq \bar{\kappa}(\xi) (\|(x, \eta) - (\bar{x}, \bar{\eta})\| + \|u - \bar{u}\|), \end{aligned} \quad (27)$$

where $\bar{\kappa}(\xi) := \max\{\kappa(\xi), \max_{(x,\eta)} \|\nabla_{(x,\eta)} g_i(\bar{x}, \bar{\eta}, \xi)\|\}$. Due to the compactness of $\mathcal{E} \setminus \mathcal{E}^\epsilon(\bar{x}, \bar{\eta})$, we claim through the finite covering theorem that there exists a unified $\delta_{(\bar{x}, \bar{\eta})} > 0$ such that (27) holds for all $(x, \eta, u) \in B((\bar{x}, \bar{\eta}, \bar{u}), \delta_{(\bar{x}, \bar{\eta})})$ and all $\xi \in \mathcal{E} \setminus \mathcal{E}^\epsilon((\bar{x}, \bar{\eta}))$. \square

Following Remark 4.1, Proposition 4.1(ii) implies that $g^o(x, \eta, \xi; u)$ is limiting almost H-calm. Indeed condition (c) of Proposition 4.1 guarantees that set \mathcal{E}_x consists a finite number of points. This is the weakest verifiable sufficient condition that we could find to ensure $\mu(\mathcal{E}_x) = 0$ and existence of a positive constant δ_x : without this condition we are unable to show almost H-calmness of $g^o(x, \eta, \xi; u)$ or limiting almost H-calmness.

Theorem 4.2 *Let $\mathcal{Z} \times \Lambda$ be a nonempty compact subset of $\mathcal{Z} \times \mathbb{R}_+$ and $H_N(x, \eta)$ be defined as in (4). Suppose, in addition to conditions of Proposition 4.1, that:* (a) $c_i(x, \xi)$ *is locally Lipschitz continuous w.r.t. x for every ξ with modulus $\kappa(\xi)$, where $\mathbb{E}[\kappa(\xi)] < \infty$, and (b) the support set of ξ is bounded. Let*

$$\begin{aligned} R_1^N(x, \eta, \lambda) &:= \mathbb{D} \left(\begin{pmatrix} \nabla f(x_N) \\ 0 \end{pmatrix} + \begin{pmatrix} 0 \\ \lambda \end{pmatrix} + \frac{\lambda}{\alpha N} \sum_{i=1}^N \partial_{(x,\eta)} g(x, \eta, \xi^i), \right. \\ &\quad \left. \begin{pmatrix} \nabla f(x_N) \\ 0 \end{pmatrix} + \mathbb{E} \left[\begin{pmatrix} 0 \\ 1 \end{pmatrix} + \frac{1}{\alpha} \partial_{(x,\eta)} g(x, \eta, \xi) \right] \lambda \right) \end{aligned}$$

and

$$R_2^N(x, \eta, \lambda) := \|\max(-H_N(x, \eta), \lambda) - \max(-H(x, \eta), \lambda)\|.$$

Then, with probability approaching one exponentially fast with the increase of sample size N , $\sup_{(x,\eta,\lambda) \in \mathcal{Z} \times \Lambda} R_1^N(x, \eta, \lambda)$ and $\sup_{(x,\eta,\lambda) \in \mathcal{Z} \times \Lambda} R_2^N(x, \eta, \lambda)$ tend to 0.

Proof It is well known that for any compact sets A, B , $\mathbb{D}(A, B) := \inf_{t>0} \{t : A \subset B + tB\}$. Using this equivalence definition, one can easily derive that, for any set C , $\mathbb{D}(A, B) \leq \mathbb{D}(A, C) + \mathbb{D}(C, B)$ and $\mathbb{D}(A + C, B + C) \leq \mathbb{D}(A, B)$. Consequently, we have

$$R_1^N(x, \eta, \lambda) \leq \mathbb{D} \left(\frac{1}{\alpha N} \sum_{i=1}^N \partial_{(x, \eta)} g(x, \eta, \xi^i) \lambda, \frac{1}{\alpha} \mathbb{E} [\partial_{(x, \eta)} g(x, \eta, \xi)] \lambda \right).$$

Since the Clarke subdifferential is convex and compact set-valued, we can use the Hörmander's formula (14) to reformulate the right-hand side of the inequality above as

$$\max_{\|u\| \leq 1} \left[\sigma \left(\frac{1}{\alpha N} \sum_{i=1}^N \partial_{(x, \eta)} g(x, \eta, \xi^i) \lambda, u \right) - \sigma \left(\frac{1}{\alpha} \mathbb{E} [\partial_{(x, \eta)} g(x, \eta, \xi)] \lambda, u \right) \right]$$

and, by virtue of the property of support function (see [27]) and [28, Proposition 3.4], further as

$$\frac{1}{\alpha} \max_{\|u\| \leq 1} \left[\frac{1}{N} \sum_{i=1}^N g_{(x, \eta)}^o(x, \eta, \xi^i; u) - \mathbb{E} [g_{(x, \eta)}^o(x, \eta, \xi; u)] \right] \lambda. \quad (28)$$

Let $\Delta_N(x, \eta, u) := \frac{1}{N} \sum_{i=1}^N g_{(x, \eta)}^o(x, \eta, \xi^i; u) - \mathbb{E} [g_{(x, \eta)}^o(x, \eta, \xi; u)]$. Since Λ is a compact set, there exists a positive number M such that $\sup_{\lambda \in \Lambda} \lambda \leq M$. Consequently,

$$\begin{aligned} & \text{Prob} \left\{ \sup_{\lambda \in \Lambda, (x, \eta) \in \mathcal{Z}, \|u\| \leq 1} \Delta_N(x, \eta, u) \lambda \geq \epsilon \right\} \\ & \leq \text{Prob} \left\{ M \sup_{(x, \eta) \in \mathcal{Z}, \|u\| \leq 1} \Delta_N(x, \eta, u) \geq \epsilon \right\}. \end{aligned} \quad (29)$$

By Proposition 4.1, $g_{(x, \eta)}^o(x, \eta, \xi; u)$ is almost H-calm with modulus $\kappa(\xi)$ and order 1 and $\mathbb{E} [g_{(x, \eta)}^o(x, \eta, \xi; u)]$ is a continuous function. Under condition (b), the moment generating function

$$M_g(x, \eta) := \mathbb{E} [e^{(g_{(x, \eta)}^o(x, \eta, \xi; u) - \mathbb{E} [g_{(x, \eta)}^o(x, \eta, \xi; u)])t}]$$

and

$$M_\kappa(t) := \mathbb{E} [e^{[\kappa(\xi) - \mathbb{E} [\kappa(\xi)]]t}]$$

are finite valued for t close to 0. By Theorem 8.1, for any $\epsilon > 0$, there exist constants $c_1(\epsilon) > 0$ and $\beta_1(\epsilon) > 0$ (independent of N) such that

$$\begin{aligned} & \text{Prob} \left\{ \sup_{(x, \eta) \in \mathcal{Z}, \|u\| \leq 1} \frac{1}{N} \sum_{i=1}^N g_{(x, \eta)}^o(x, \eta, \xi^i; u) - \mathbb{E} [g_{(x, \eta)}^o(x, \eta, \xi; u)] \geq \frac{\epsilon \alpha}{M} \right\} \\ & \leq c_1(\epsilon) e^{-N\beta_1(\epsilon)}. \end{aligned}$$

Combining the inequality above with (28) and (29), we have

$$\begin{aligned} & \text{Prob} \left\{ \sup_{(x, \eta, \lambda) \in \mathcal{Z} \times \Lambda} R_1^N(x, \eta, \lambda) \geq \epsilon \right\} \\ & \leq \text{Prob} \left\{ \sup_{(x, \eta) \in \mathcal{Z}, \|u\| \leq 1} \frac{1}{N} \sum_{i=1}^N g(x, \eta, \xi^i, u) - \mathbb{E}[g(x, \eta, \xi, u)] \geq \frac{\epsilon \alpha}{M} \right\} \\ & \leq c_1(\epsilon) e^{-\beta_1(\epsilon)N}. \end{aligned} \quad (30)$$

On the other hand, since $R_2^N(x, \eta, \lambda) \leq \|H(x, \eta) - H_N(x, \eta)\|$, under conditions (a) and (b), it follows by [12, Theorem 3.1] that, for any small positive number $\epsilon > 0$, there exist positive constants $c_2(\epsilon)$ and $\beta_2(\epsilon)$ (independent of N) such that

$$\begin{aligned} \text{Prob} \left\{ \sup_{(x, \eta, \lambda) \in \mathcal{Z} \times \Lambda} R_2^N(x, \eta, \lambda) \geq \epsilon \right\} & \leq \text{Prob} \left\{ \sup_{(x, \eta) \in \mathcal{Z}} \|H(x, \eta) - H_N(x, \eta)\| \geq \epsilon \right\} \\ & \leq c_2(\epsilon) e^{-\beta_2(\epsilon)N}. \end{aligned} \quad (31)$$

The proof is complete. \square

Let

$$\Gamma(x, \eta, \lambda) := \begin{pmatrix} \begin{pmatrix} \nabla f(x) \\ 0 \end{pmatrix} + \mathbb{E} \left[\begin{pmatrix} 0 \\ 1 \end{pmatrix} + \frac{1}{\alpha} \partial_{(x, \eta)} g(x, \eta, \xi) \right] \lambda \\ \max(-H(x, \eta), \lambda) \end{pmatrix}$$

and

$$\mathcal{G}(x, \eta) := \begin{pmatrix} \mathcal{N}_{\mathcal{X} \times \mathbb{R}}(x, \eta) \\ 0 \end{pmatrix}.$$

We can rewrite (9) as a stochastic generalized equation

$$0 \in \Gamma(x, \eta, \lambda) + \mathcal{G}(x, \eta).$$

Likewise, we can rewrite the KKT conditions (12) of the SAA problem as follows:

$$0 \in \Gamma_N(x, \eta, \lambda) + \mathcal{G}(x, \eta),$$

where

$$\Gamma_N(x, \eta, \lambda) := \begin{pmatrix} \begin{pmatrix} \nabla f(x) \\ 0 \end{pmatrix} + \left(\begin{pmatrix} 0 \\ 1 \end{pmatrix} + \frac{1}{\alpha N} \sum_{i=1}^N \partial_{(x, \eta)} g(x, \eta, \xi^i) \right) \lambda \\ \max(-H_N(x, \eta), \lambda) \end{pmatrix}.$$

It is easy to see that Theorem 4.2 implies

$$\sup_{(x, \eta, \lambda) \in \mathcal{Z} \times \Lambda} \mathbb{D}(\Gamma_N(x, \eta, \lambda), \Gamma(x, \eta, \lambda)) \rightarrow 0 \quad (32)$$

a.s. $N \rightarrow \infty$.

Theorem 4.3 Assume the settings and conditions of Theorem 4.2. Under Assumption 4.1, for any $\epsilon > 0$, there exist positive constants $c(\epsilon)$ and $\beta(\epsilon)$ independent of N such that

$$\text{Prob}\{d((x_N, \eta_N), Z^*) \geq \epsilon\} \leq c(\epsilon)e^{-N\beta(\epsilon)},$$

where (x_N, η_N) denotes the KKT point satisfying (12) and Z^* denotes the set of weak Clarke stationary points characterized by (9).

Proof The thrust of the proof is to use (32) and [12, Lemma 4.2]. To this end, we need to verify the outer semicontinuity of $\Gamma(x, \eta, \lambda)$. Observe that $\Gamma(x, \eta, \lambda)$ consists of two parts:

$$\begin{pmatrix} \nabla f(x) \\ 0 \end{pmatrix} + \mathbb{E} \left[\begin{pmatrix} 0 \\ 1 \end{pmatrix} + \frac{1}{\alpha} \partial_{(x, \eta)} g(x, \eta, \xi) \right] \lambda$$

and $\max(H(x, \eta), \lambda)$. Since f is continuously differentiable and $\partial_{(x, \eta)} g(x, \eta, \xi)$ is outer semicontinuous w.r.t. x, η for almost every ξ and integrably bounded, it follows by Aumann [17, Corollary 5.2],

$$\begin{pmatrix} \nabla f(x) \\ 0 \end{pmatrix} + \mathbb{E} \left[\begin{pmatrix} 0 \\ 1 \end{pmatrix} + \frac{1}{\alpha} \partial_{(x, \eta)} g(x, \eta, \xi) \right] \lambda$$

is outer semicontinuous. Moreover, since $H(x, \eta)$ is a continuous function, $\max(-H(x, \eta), \lambda)$ is continuous w.r.t. (x, η) , and λ . Therefore, $\Gamma(x, \eta, \lambda)$ is outer semicontinuous. The rest follows from (32) and [12, Lemma 4.2]. \square

It is important to note that the constants $c(\epsilon)$ and $\beta(\epsilon)$ in Theorem 4.3 may be significantly different from their counterparts in Theorem 4.2. To establish a precise relationship of these constants, we will need more information about the sensitivity of the true problem at the stationary points. One possibility is to look into the metric regularity type condition for set-valued mapping $\Gamma(x, \eta, \lambda) + \mathcal{G}(x, \eta)$. If there exists positive constants C and γ such that

$$d(x, \eta, Z^*) \leq Cd(0, \Gamma(x, \eta, \lambda) + \mathcal{G}(x, \eta))^\gamma$$

for x, η close to Z^* , then we can establish

$$d(x, \eta, Z^*) \leq C(\|R_1^N(x, \eta, \lambda)\| + \|R_2^N(x, \eta, \lambda)\|)^\gamma.$$

We refer interested readers to [29] for recent discussions on metric regularity. Under this circumstance, the constants $c(\epsilon)$ and $\beta(\epsilon)$ can be easily expressed in terms of $c(\epsilon) := c_1(\epsilon) + c_2(\epsilon)$ and $\beta(\epsilon) := \min(\beta_1(\epsilon), \beta_2(\epsilon))$ in Theorem 4.2. Moreover, following [13, Remark 3.1], under some additional conditions on the moment functions, we can obtain an estimation of sample size through (30) and (31), that is there exists a constant $\sigma > 0$ such that for any $\epsilon > 0$, $\text{Prob}\{d((x_N, \eta_N), Z^*) \geq \epsilon\} \leq \beta$ holds when

$$N \geq \frac{O(1)\sigma^2}{\epsilon^2} \left[n \ln \left(O(1)D \left(\frac{4\mathbb{E}[\kappa(\xi)]}{\epsilon} \right) \right)^{\frac{1}{\gamma}} + \ln \left(\frac{1}{\beta} \right) \right],$$

where $D := \sup_{(x', \eta'), (x, \eta) \in \mathcal{Z}} \|(x', \eta') - (x, \eta)\|$ is the diameter of \mathcal{Z} and $O(1)$ is a generic constant. We leave this to interested readers as it involves complex technical details.

5 DC-Approximation

Although CVaR approximation is known to be the “best” convex approximation method of chance constraints, as commented in [4], it is a convex conservative approximation, which means that there exists a gap between the CVaR approximation and the true constraint. In [14, 30], Hong et al. proposed a DC-approximation method for a joint chance constraint. The numerical tests show that the DC-approximation scheme displays better results than CVaR approximation scheme. Hong et al. also showed almost sure convergence of optimal solution of subproblems in their algorithm called *sequential convex approximations* (SCA). It is easy to extend their discussion to convergence of stationary points. In this section, we only list the similar convergence results as above section under the DC-approximation scheme but omit the proofs for the limitation of the length of the paper. We refer interested readers to [31, Sect. 4] for details.

The formulation of the DC-approximation problem is defined as follows:

$$\min_{x \in \mathcal{X}} f(x) \quad \text{s.t.} \quad \inf_{t > 0} \frac{1}{t} (\mathbb{E}[p(c(x, \xi) + t)] - \mathbb{E}[p(c(x, \xi))]) \leq \alpha, \quad (33)$$

where t is a positive number. In [14], Hong et al. use ε -approximation of problem (33) by setting $t = \varepsilon$. The formulation of ε -approximation problem is

$$\min_{x \in \mathcal{X}} f(x) \quad \text{s.t.} \quad H_\varepsilon^{\text{DC}}(x) := \mathbb{E}[p(c(x, \xi) + \varepsilon)] - \mathbb{E}[p(c(x, \xi))] - \varepsilon\alpha \leq 0, \quad (34)$$

and its SAA problem as

$$\min_{x \in \mathcal{X}} f(x) \quad \text{s.t.} \quad H_N^{\text{DC}}(x) := \frac{1}{N} \sum_{j=1}^N (p(c(x, \xi^j) + \varepsilon) - p(c(x, \xi^j))) - \varepsilon\alpha \leq 0. \quad (35)$$

Hong et al. also prove that when $\varepsilon \downarrow 0$, the KKT point of problem (34) converges to that of problem (1).

Let $\lambda \geq 0$ and define the Lagrange function of problem (33):

$$\mathcal{L}^{\text{DC}}(x, \lambda) := f(x) + \lambda H_\varepsilon^{\text{DC}}(x).$$

In order to derive the KKT conditions of problem (33), we need some assumptions:

Assumption 5.1 $c_i(x, \xi)$, $i = 1, \dots, m$, is locally Lipschitz continuous w.r.t. x with modulus $\kappa_i(\xi)$ where $\mathbb{E}[\kappa_i(\xi)] < \infty$.

Assumption 5.2 $c(\cdot, \xi)$ is differentiable on \mathcal{X} for a.e. ξ .

Assumption 5.3 Let $F(t, x) := \text{Prob}\{c(x, \xi) \leq t\}$. There exists a constant $\varrho > 0$ such that $F(t, x)$ is continuously differentiable on $(-\varrho, \varrho) \times \mathcal{X}$.

Assumptions 5.1–5.3 are used and discussed in [14].

Definition 5.1 Problem (34) is said to satisfy the *differential constraint qualification* at a feasible point x iff there exists $d \in \mathbb{R}^n$ such that $\nabla_x H_\varepsilon^{\text{DC}}(x)^T d < -\delta$, where $\delta > 0$ is a positive constant.

Proposition 5.1 Let $x^* \in \mathcal{X}$ be a local optimal solution to the true problem (34). Let Assumptions 2.1, 5.1–5.3 hold and the differential constraint qualification be satisfied at x^* . Then there exists a $\lambda^* \in R_+$ such that

$$\begin{cases} 0 \in \nabla f(x^*) + \lambda^* \nabla_x H_\varepsilon^{\text{DC}}(x^*) + \mathcal{N}_{\mathcal{X}}(x^*), \\ 0 \leq -H_\varepsilon^{\text{DC}}(x^*) \perp \lambda^* \geq 0. \end{cases} \quad (36)$$

Proposition 5.2 Let $x_N \in \mathcal{X}$ be a local optimal solution to the sample average approximation problem (35). Let \hat{X} denote a subset of \mathcal{X} such that

$$\lim_{N \rightarrow \infty} d(x_N, \hat{X}) \rightarrow 0, \quad \text{w.p.1.}$$

Suppose that Assumptions 2.1, 5.1–5.3 hold, \hat{X} is bounded and the differential constraint qualification holds at every point $x \in \hat{X}$. Then w.p.1 problem (35) satisfies the differential constraint qualification for N sufficiently large, and there exists $\lambda_N \in R_+$ such that

$$\begin{cases} 0 \in \nabla f(x_N) + \Phi_N(x_N) \lambda_N + \mathcal{N}_{\mathcal{X}}(x_N), \\ 0 \leq -H_N^{\text{DC}}(x_N) \perp \lambda_N \geq 0, \end{cases} \quad (37)$$

where $\Phi_N(x) := \frac{1}{N} \sum_{j=1}^N (\nabla_x c_{i(x)}(x, \xi^j) \cdot \mathbb{1}_{(-\varepsilon, +\infty)}(c(x, \xi^j)) - \nabla_x c_{i(x)}(x, \xi^j) \cdot \mathbb{1}_{(0, +\infty)}(c(x, \xi^j)))$.

We call a tuple (x^*, λ^*) satisfying (36) a *KKT pair* of problem (34), x^* a *stationary point* and λ^* the corresponding Lagrange multiplier and a tuple (x_N, λ_N) satisfying (37) a *KKT pair* of problem (35). We also note that such a (x^*, λ^*) can be thought as an ε -KKT pair of problem (33).

We make a blanket assumption that throughout the rest of this section the conditions of Proposition 5.2 hold.

Assumption 5.4 There exists a compact subset $X \times \Lambda \subset \mathcal{X} \times R_+$ and a positive number N_0 such that w.p.1 problem (35) has a KKT pair $(x_N, \lambda_N) \in X \times \Lambda$ for $N \geq N_0$.

Theorem 5.1 Let $\{(x_N, \lambda_N)\}$ be a sequence of KKT pairs of problem (35) and (x^*, λ^*) be a cluster point. Suppose that Assumptions 5.1–5.4 hold. Then w.p.1 (x^*, λ^*) is a KKT pair of the true problem (34), which satisfies the KKT system (36).

Theorem 5.2 Let $X \times \Lambda$ be a nonempty compact subset of $\mathcal{X} \times \mathbb{R}_+$. Let $H_N^{\text{DC}}(x, \eta)$ be defined as in (35) and $\Phi_N(x)$ be defined as in (37). Assume: (a) Assumptions 5.1, 5.3, and 5.4 hold; (b) the support set of ξ is bounded; (c) $c_i(x, \xi)$, $i = 1, \dots, m$, is continuously differentiable w.r.t. (x, ξ) and twice continuously differentiable w.r.t. x for almost every $\xi \in \Xi$; (d) there exists an integrable function $\kappa : \Xi \rightarrow \mathbb{R}$ such that $\nabla_x c_i(\cdot, \xi)$ is locally Lipschitz continuous with modulus $\kappa(\xi)$ for every $\xi \in \Xi$ where $\mathbb{E}[\kappa(\xi)] < \infty$; (e) $\hat{\Xi}(x)$, $\hat{\Xi}_{i,j}(x)$, $\tilde{\Xi}(x)$ and $\tilde{\Xi}_{i,j}(x)$ are compact,

$$\nabla_{\xi}(\hat{c}_i(x, \xi) - \hat{c}_j(x, \xi)) \neq 0, \quad \forall \xi \in \hat{\Xi}_{i,j}(x),$$

and

$$\nabla_{\xi}(c_i(x, \xi) - c_j(x, \xi)) \neq 0, \quad \forall \xi \in \tilde{\Xi}_{i,j}(x)$$

hold for all $i, j \in \{1, \dots, m+1\}$, $i \neq j$. Then, for any $\epsilon > 0$, there exist positive constants $C(\epsilon)$ and $\beta(\epsilon)$ independent of N such that

$$\text{Prob}\{d(x_N, X^*) \geq \epsilon\} \leq C(\epsilon)e^{-N\beta(\epsilon)},$$

where x_N denotes the Clarke stationary points characterized by (37) and X^* denotes the set of Clarke stationary points characterized by (36).

6 Numerical Tests

We have carried out a number of numerical experiments on the approximation scheme for (1) in Matlab 7.9.0 installed in a PC with Windows XP operating system. To deal with joint chance constraint, we apply CVaR method to approximate it and then reformulate the latter as (3). In the tests, we apply SAA method to problem (3) and employ the random number generator *rand* in Matlab 7.9.0 to generate the samples and solver *fmincon* to solve the SAA problem (4). Since our focus in this paper is on convergence analysis, then our numerical test is to show the convergence behavior with the sample size increase.

Let $x = (x_1, \dots, x_d)^T$ denote a d -dimensional vector in \mathbb{R}^d and $\xi \in \mathbb{R}^{d \times m}$ be a matrix of random variables. We consider the following norm optimization problem:

$$\begin{aligned} \min \quad & -\sum_{j=1}^d x_j \\ \text{s.t.} \quad & \text{Prob}\left\{\sum_{j=1}^d \xi_{ij}^2 x_j^2 \leq 100, \quad i = 1, \dots, m\right\} \geq 1 - \alpha, \\ & x_j \geq 0, \quad j = 1, \dots, d. \end{aligned} \quad (38)$$

Note that (38) is a joint chance constraint problem; see [14, Sect. 5.1] for details.

Let $c_i(x, \xi) := \sum_{j=1}^d \xi_{ij}^2 x_j^2 - 100$, for $i = 1, \dots, m$. For any $x \neq 0$, $c_i(x, \xi)$ is a continuous random variable and $c_i(x, \xi) = c_j(x, \xi)$ with probability 0.

We consider the case when $\xi_{ij}, i = 1, \dots, m, j = 1, \dots, d$, are independent and identically distributed random variables with stand normal distribution. Via a similar argument to that in [14], we can work out the optimal solution of (38) $x^* = \{x_1^*, \dots, x_d^*\}$, where

$$x_i^* = \frac{10}{F_{\chi_d^2}^{-1}((1 - \alpha)^{1/m})}, \quad i = 1, \dots, d,$$

and $F_{\chi_d^2}^{-1}$ denotes the inverse distribution function of a chi-square distribution with d degrees of freedom. We use x^* as a benchmark in the CVaR schemes.

For $d = 2$ and $m = 2$, the true optimal solution of problem (38) is $x^* = (4.10, 4.10)$ with optimal value $f^* = -8.20$. We set $\alpha = 0.1$ and $\epsilon = 0.05^2$ and perform comparative analysis with respect to the sample size from 100 to 2800 with increment 300. For every fixed sample size, 50 independent tests are carried out each of which solves the SAA problem and yields an approximation solution. We solve the problem (38) under the two approximation schemes, CVaR and DC, and use the numerical solution of CVaR problem, x_{CVaR} , as a start point of DC problem. We carry out those numerical experiments on this problem with gradient-based Monte Carlo method and Algorithm SCA; see [14], in Matlab 7.9.0 installed in a PC with windows XP where CVaR problem and every convex subproblem generated by Algorithm SCA are solved by *fmincon*.

We use a vertical interval to indicate the range of the 50 approximate optimal values and optimal solutions based on the CVaR and DC-approximation schemes. As sample size increases, we observe a trend of convergence of the range of the optimal values and solutions. Figure 1 depicts the convergence of optimal values due to CVaR and DC. Figure 2 displays the convergence of two components of the approximate optimal solution based on CVaR approximation while Fig. 3 shows the similar convergence of two components of the approximate optimal solution based on DC approximation. All of those figures show that when sample size increases from 100 to 2800, both optimal values and optimal solutions converge quickly and when sample size reaches 1300 there is not substantial changes of these quantities. Although the numerical tests show that the DC approximation methods approach the joint chance constrained problem better, the CPU time of solving DC approximation scheme via Algorithm SCA is much more than its CVaR counterpart.

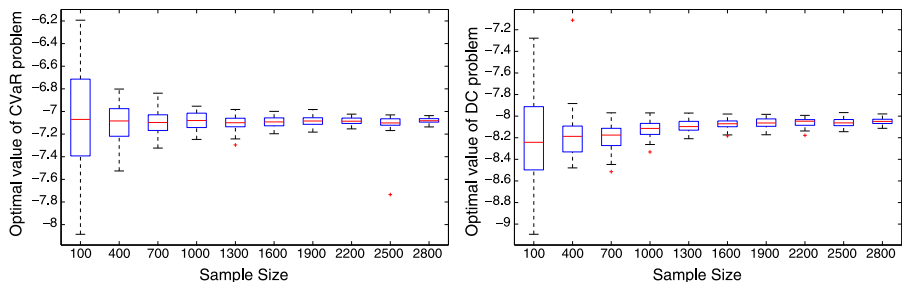


Fig. 1 The convergence of the optimal values of CVaR and DC-approximation

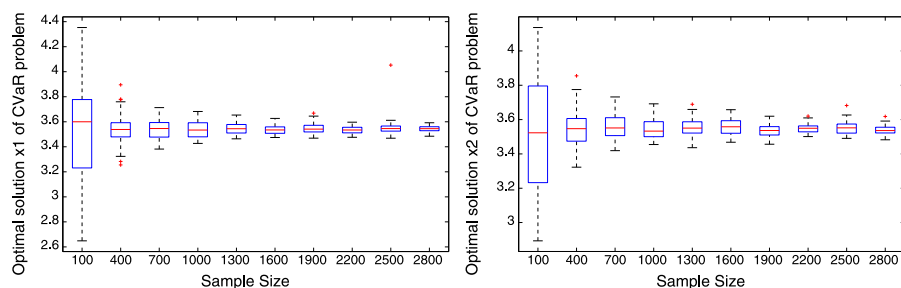


Fig. 2 The convergence of the optimal solutions of CVaR approximation

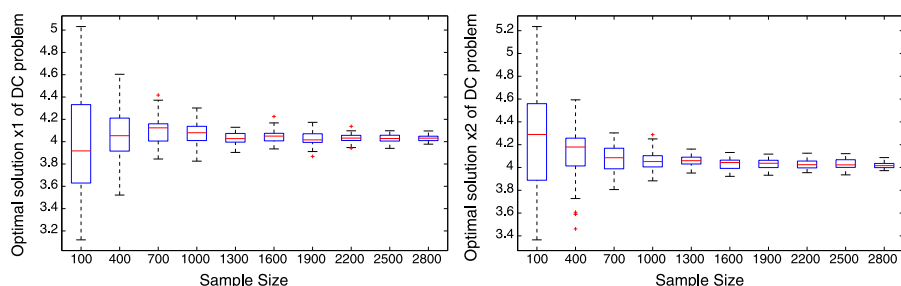


Fig. 3 The convergence of the optimal solutions of DC-approximation

7 Concluding Remarks

In this paper, we present a detailed asymptotic analysis of stationary points of a stochastic program with a specific CVaR constraint which approximates a joint chance constraint. CVaR has been widely used as a risk measure in finance, engineering, and management sciences [32–34] and it is of immense interest to investigate optimization problem with CVaR as a constraint rather than an objective as in the literature [3, 7, 8]. Our analysis has addressed the needed asymptotic consistency and numerical efficiency/tractability of estimators of optimal solutions and stationary points under some moderate conditions when the well-known Monte Carlo method is applied to solve the CVaR constrained problem. The analysis is extended to the case when the joint chance constraint is approximated via DC-programming although the details are left in an earlier version [31] due to the limitation of length of the paper.

Note that the established convergence results are problem specific, that is, the CVaR or DC functions are constructed from the underlying functions of the joint chance constraints. However, we envisage that this kind of convergence may be established when SAA is applied to a stochastic problem with a general CVaR or DC constraint. For instance, in the analysis of exponential convergence, we employ the concept of almost H-calmness which is satisfied by a broad class of piecewise smooth random functions. We leave this to interested readers.

Note also that joint chance constrained problem is intrinsically nonconvex, therefore, as a convex approximation, CVaR does not necessarily give rise to the best approximation to the true problem; indeed, it has already been shown that DC-approximation delivers a better approximation [14] although the former is much easier to solve due to the convex constraint.

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Appendix

In this [Appendix](#), we strengthen [13, Theorem 3.1] by weakening a boundedness condition imposed on the random function.

Theorem 8.1 *Let $\phi : \mathbb{R}^n \times \mathcal{E} \rightarrow \mathbb{R}$ be a real valued lower semicontinuous function, $\xi : \Omega \rightarrow \mathcal{E} \subset \mathbb{R}^k$ a random vector defined on probability space (Ω, \mathcal{F}, P) and $\psi(x) := \mathbb{E}[\phi(x, \xi)]$. Let $\mathcal{X} \subset \mathbb{R}^n$ be a compact subset of \mathbb{R}^n . Assume: (a) for every $x \in \mathcal{X}$ the moment generating function $M_x(t) := \mathbb{E}\{e^{t[\phi(x, \xi) - \psi(x)]}\}$ is finite valued for all t in a neighborhood of zero. (b) $\psi(x)$ is continuous on \mathcal{X} , (c) $\phi(x, \xi)$ is bounded by an integrable function $L(\xi)$ and the moment generating function $\mathbb{E}[e^{(L(\xi) - \mathbb{E}[L(\xi)])t}]$ is finite valued for t close to 0. Then the following statements hold.*

- (i) *If $\phi(\cdot, \xi)$ is almost H -clam from above at every point $x \in \mathcal{X}$ with modulus $\kappa_x(\xi)$ and order γ_x , and the moment generating function $\mathbb{E}[e^{\kappa_x(\xi)t}]$ is finite valued for t close to 0, then for every $\epsilon > 0$, there exist positive constants $c(\epsilon)$ and $\beta(\epsilon)$, independent of N , such that*

$$\text{Prob}\left\{\sup_{x \in \mathcal{X}} (\psi_N(x) - \psi(x)) \geq \epsilon\right\} \leq c(\epsilon)e^{-N\beta(\epsilon)}. \quad (39)$$

- (ii) *If $\phi(\cdot, \xi)$ is almost H -clam from below at every point $x \in \mathcal{X}$ with modulus $\kappa_x(\xi)$ and order γ_x , and the moment generating function $\mathbb{E}[e^{\kappa_x(\xi)t}]$ is finite valued for t close to 0, then for every $\epsilon > 0$, there exist positive constants $c(\epsilon)$ and $\beta(\epsilon)$, independent of N , such that*

$$\text{Prob}\left\{\sup_{x \in \mathcal{X}} (\psi_N(x) - \psi(x)) \leq -\epsilon\right\} \leq c(\epsilon)e^{-N\beta(\epsilon)}. \quad (40)$$

- (iii) *If $\phi(\cdot, \xi)$ is almost H -clam at every point $x \in \mathcal{X}$ with modulus $\kappa_x(\xi)$ and order γ_x , and the moment generating function $\mathbb{E}[e^{\kappa_x(\xi)t}]$ is finite valued for t close to 0, then for every $\epsilon > 0$, there exist positive constants $c(\epsilon)$ and $\beta(\epsilon)$, independent of N , such that*

$$\text{Prob}\left\{\sup_{x \in \mathcal{X}} |\psi_N(x) - \psi(x)| \geq \epsilon\right\} \leq c(\epsilon)e^{-N\beta(\epsilon)}. \quad (41)$$

Due to the limitation of the length of the paper, we omit the proof which can be found in [31].

Note that the exponential convergence is derived for the case when ξ satisfies a continuous distribution. In the case when ξ satisfies a discrete distribution, the concept of almost H-calmness is no longer applicable. However, the uniform exponential convergence may be established in an entirely different way for a class of random function which is uniformly bounded over a considered compact set. We leave this to interested readers.

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